

On an ambiguity in the concept of partial and total derivatives in classical analysis

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(February 1, 2008)

Abstract

Ambiguity is shown in the context of the differential calculus of several variables and with the help of the language of category theory, a way to solve it in its most general form is offered. It is also shown that this new definition is related to other well-known definitions in the literature.

I. INTRODUCTION

The difference between the functions:

$$E[x_1(t), \dots, x_{n-1}(t), t] = \text{def } E[\mathbf{r}(t), t], \quad E(x_1, \dots, x_{n-1}, t) = \text{def } E(\mathbf{r}, t)$$

is usually not remarked in the literature, and for this reason we can often write down meaningless symbols like:

$$\frac{\partial}{\partial t} E[\mathbf{r}(t), t], \tag{1}$$

and

$$\frac{d}{dt} E(\mathbf{r}, t). \tag{2}$$

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Ambiguities in the “notation” for partial differentiation has been remarked by Arnold [1] p. 226 (p. 258 in English translation) without further development. The symbols (1), (2) are meaningless, because the process denoted by the operator of *partial* differentiation can be applied only to functions of several *independent* variables and $E[\mathbf{r}(t), t]$ is not *such* a function. Meanwhile, the operator of *total* differentiation with respect to given variable can be formally applied to functions of one variable only. However, we have a well-known formula to relate both concepts:

$$\frac{d}{dt}E = (\mathbf{V} \cdot \nabla)E + \frac{\partial}{\partial t}E \quad (3)$$

(here $\mathbf{V} = \frac{d\mathbf{r}}{dt}$).

Let us show that, in this form, Eq. (3) cannot be correct. What is the correct argument for the symbol E in both sides? If we say that the correct argument for both sides is $[\mathbf{r}(t), t]$ we get the chain of symbols (1), but in this case, the operator of a partial differentiation would indicate that we must construct a new function in the form $(\partial E / \partial t)$, hence we use the following procedure:

$$\lim_{\Delta t \rightarrow 0} \left\{ \frac{E \left[\mathbf{r}(t) + \Delta t \frac{d\mathbf{r}(t)}{dt}, t + \Delta t \right] - E[\mathbf{r}(t), t]}{\Delta t} \right\}. \quad (4)$$

But this is the definition of total differentiation! Thus, the symbols of total and of partial differentiation denote the same process, therefore, because E is the same function on both sides of the equation, we get:

$$(\mathbf{V} \cdot \nabla)E[\mathbf{r}(t), t] = 0 \quad (5)$$

always. But even if the procedure which we followed were correct (which it is not, of course!), this equation is not correct for E as a function of the functions $\mathbf{r}(t)$, because the partial differentiation would involve increments of the functions $\mathbf{r}(t)$ in the form $\mathbf{r}(t) + \Delta\mathbf{r}(t)$ and we do not know how we must interpret this increment because we have two options: *either* $\Delta\mathbf{r}(t) = \mathbf{r}(t) - \mathbf{r}^*(t)$, *or* $\Delta\mathbf{r}(t) = \mathbf{r}(t) - \mathbf{r}(t^*)$. Both are different processes because the first one involves changes in the functional form of the functions $\mathbf{r}(t)$, while the second involves changes in the position along the path defined by $\mathbf{r} = \mathbf{r}(t)$ but preserving the same functional form. Hence, it is clear that we have here different concepts. If we remember the definition of partial differentiation, we can see where the mistake is: “*the*

symbol: $\frac{\partial}{\partial t} E(\mathbf{r}, t)$ means that we take the variations of t when the values of \mathbf{r} are constant". It means that we make the only change $t + \Delta t$ in the function. But this is only possible if the coordinates \mathbf{r} are independent from t . Hence, we can see that the correct argument cannot be $[\mathbf{r}(t), t]$, because, as we have shown, this supposition leads to the incorrect result (5). If we make the other supposition, that the correct argument is (\mathbf{r}, t) we can get the same conclusion, i.e., equation (5). Hence, *none of these suppositions is correct*. What is the solution, then? Actually, in the equation (3) we have two *different* functions: on the left hand side we have the function $E[\mathbf{r}(t), t]$ defined on a *curve* in a n -surface and on the right hand side we have the function $E(\mathbf{r}, t)$ defined on the *all* n -surface, which obviously are *quite* different functions, while we have a limiting procedure to get a unification of concepts in the realm of functions of one variable.

Now let us introduce the following notation:

$$f = E \circ \mathbf{p}, \quad (6)$$

where the symbol "o" means a composition of functions and where

$$E : R^n \rightarrow R, \quad \mathbf{p} : R \rightarrow R^n, \quad f : R \rightarrow R.$$

It is clear that $\mathbf{p} = \mathbf{p}(t) = \{x_1(t), \dots, x_{n-1}(t), t\} = \{\mathbf{r}(t), t\}$ is a curve which lies on the n -surface where the function E is defined.

Hence we can write down the equation:

$$\frac{d}{dt} f = \lim_{x_i \rightarrow x_i(t)} \left\{ (\mathbf{V} \cdot \nabla) E + \frac{\partial E}{\partial t} \right\}$$

which shows our point more clearly: *the functions in both sides (f and E) are different functions*. Of course, we suppose that the components of the vector \mathbf{V} tend to derivatives $\frac{dx_i}{dt}$ in the limit. But here is where our grammatical distinction appears: the right hand side is evaluated in all points along the curve $\mathbf{p}(t)$, that is:

$$(\mathbf{V} \cdot \nabla) E \Big|_{x_i=x_i(t)} + \frac{\partial E}{\partial t} \Big|_{x_i=x_i(t)}.$$

Let us explain the distinction as follows: the operator of the total differentiation is just a differentiation of a function which can depend on one independent variable, and the operator of the partial differentiation is just a partial differentiation of a function which

can depend on several independent variables. An obvious question immediately arises: *what is the relation between these domains?* Obviously, the function of one variable is one entity and the function of several variables is a different one. The relation lies in the evaluation of the function obtained by partial differentiation in points *along* the curve $\mathbf{p}(t)$. Or in more general terms, we must have the validity of the following condition: for all $\varepsilon > 0$ there is a $\delta[\varepsilon, \mathbf{p}(t)] > 0$ such that if we take a point in the ball:

$$|\hat{\mathbf{r}} - \mathbf{p}(t)| < \delta[\varepsilon, \mathbf{p}(t)],$$

where $\hat{\mathbf{r}} = \{x_1, \dots, x_{n-1}, t\} = \{\mathbf{r}, t\}$, then

$$\left| \left[\left\{ (\mathbf{V} \cdot \nabla) E \right\}(\mathbf{r}, t) + \left\{ \frac{\partial E}{\partial t} \right\}(\mathbf{r}, t) \right] - \left[\left\{ (\mathbf{V} \cdot \nabla) E \right\}[\mathbf{p}(t)] + \left\{ \frac{\partial E}{\partial t} \right\}[\mathbf{p}(t)] \right] \right| < \varepsilon,$$

where, of course, $E = E(\hat{\mathbf{r}}) = E(\mathbf{r}, t)$.

We have not supposed, of course, that we have an uniform continuity. The abbreviated form of this condition is:

$$\frac{d}{dt} f(t) = \lim_{\hat{\mathbf{r}} \rightarrow \mathbf{p}(t)} \left\{ (\mathbf{V} \cdot \nabla) E(\mathbf{r}, t) + \frac{\partial}{\partial t} E(\mathbf{r}, t) \right\}. \quad (7)$$

The distinction between¹ $E(\mathbf{r}, t)$ and $E[\mathbf{r}(t), t]$ is important in some physical contexts, as it is shown in [2] (see, especially, Eq. (28)). The grammatical distinction is that the realm of functions of one independent variable is not the same as the realm of functions of several independent variables, and that the relation between these two realms is given by a limitation procedure.

II. SOME REMARKS RELATED TO THE FUNCTIONAL EQUATION (6)

What conditions must the relation $f = E \circ \mathbf{p}$ satisfy to make sense? It is obviously that all its elements f , E , and p have to exist, and we, in fact, must write down the more general relation:

$$f(t) = \lim_{\hat{\mathbf{r}} \rightarrow \mathbf{p}} E(\mathbf{r}, t). \quad (8)$$

¹or, that is the same, between $E(\hat{\mathbf{r}})$ and $E[\mathbf{p}(t)]$

It means that the function E must be continuous in all points of the curve \mathbf{p} .

We have to consider seven cases²:

1. Two functional form E and \mathbf{p} are known: This is the $\{E, \mathbf{p}\}$ -case;
2. E and f are known: $\{E, f\}$ -case;
3. \mathbf{p} and f are known: $\{\mathbf{p}, f\}$ -case;
4. Only E is known: $\{E\}$ -case;
5. Only f is known: $\{f\}$ -case;
6. Only \mathbf{p} is known: $\{\mathbf{p}\}$ -case;
7. All function are unknown: $\{\}$ -case.

In the $\{E, \mathbf{p}\}$ -, $\{E, f\}$ - and $\{\mathbf{p}, f\}$ -cases we can define one of the functions in terms of the other two functions, for example, in $\{E, f\}$ -case we define \mathbf{p} etc. In $\{E\}$ -, $\{f\}$ - and $\{\mathbf{p}\}$ -cases one can show that it is possible to define the other functions under certain conditions. Let us make a brief review of these classes.

$\{E\}$ -case: In this case, we only know the form of E , and we need to define the forms of the other two functions. We suppose that: $E \in C^1(R^n, R)$, $\mathbf{p} \in C^1(R, R^n)$, $f \in C^1(R, R)$. Now we write down our defining equation in the form

$$\frac{df}{dt} = \lim_{\hat{\mathbf{r}} \rightarrow \mathbf{p}} \left\{ \sum_{i=1}^{n-1} V_i(\hat{\mathbf{r}}) \frac{\partial E}{\partial x_i} + \frac{\partial E}{\partial t} \right\}, \quad (9)$$

and we propose the following two equations:

$$(a) \frac{df}{dt} = \lim_{\hat{\mathbf{r}} \rightarrow \mathbf{p}} \frac{\partial E}{\partial t} \quad \text{and} \quad (b) V_i(\hat{\mathbf{r}}) = \sum_{j=1}^{n-1} b_{ij} \frac{\partial E}{\partial x_j}, \quad (10)$$

where b_{ij} is a skew-symmetrical matrix ($b_{ij} = -b_{ji}$). This proposition has the following motivation: we define the components of the vector field \mathbf{V} by (10b) then, when we put this equation in (9), the first term on the right hand side vanishes and we get the equation (10a). This is not yet enough. We construct

²When we know all three functions, we must only check that the relation (7) is valid. This is trivial.

the curve as an integral curve of the vector field with the components (10b), i.e., the solution of the following set of equations (a non-autonomous system of differential equations):

$$\frac{dx_i}{dt} = \sum_{j=1}^{n-1} b_{ij} \frac{\partial E}{\partial x_j}. \quad (11)$$

Now with the solution of the equation (11) we have an explicit form of the curve \mathbf{p} . And we know E , hence we know its partial derivatives. Then for the function f we can write down:

$$f = \int \left(\lim_{\hat{\mathbf{r}} \rightarrow \mathbf{p}} \frac{\partial E}{\partial t} \right) dt + const. \quad (12)$$

So, with just the form of E we can define the form of the other two functions.

{p}-case: We just know the form of the curve. However, for this case we require the following conditions: $E \in C^1(R^n, R)$, $\mathbf{p} \in C^2(R, R^n)$, $f \in C^1(R, R)$. We shall follow the same methodology used in {E}-case. We know the explicit form of the curve \mathbf{p} , hence we know its derivatives in an explicit way. We use here a symbol $k_i(t) = dx_i/dt$ to denote these explicit functions. We have the following two equations from the defining relation:

$$(a) \frac{df}{dt} = \lim_{\hat{\mathbf{r}} \rightarrow \mathbf{p}} \frac{\partial E}{\partial t} \quad \text{and} \quad (b) \frac{\partial E}{\partial x_i} = \sum_{j=1}^{n-1} b_{ij} k_i(t), \quad (13)$$

where $b_{ij} = const$ for all i, j . In this case we have supposed that the components of the vector field, in the limit, are equal to the functions $k_i(t)$. The solution to these equations is:

$$E = \sum_{i,j}^{n-1} b_{ij} k_i(t) x_i + T(t) \quad \text{and} \quad f = \int \lim_{\hat{\mathbf{r}} \rightarrow \mathbf{p}} \left\{ \sum_{i,j}^{n-1} b_{ij} \frac{dk_j}{dt} x_j \right\} dt + \int \frac{dT}{dt} dt,$$

where T is an arbitrary function. In this case we have solved, first, the equation (13b) and its solution E is used to calculate the partial derivative with respect to t . Then we have calculated the limit to get the integrand to calculate f . Again, with just one entity, the curve, we can define the other two functions in the functional equation (6).

{f}-case: We just know the form of f . The defining relation is written as:

$$H(t) = \lim_{\hat{\mathbf{r}} \rightarrow \mathbf{p}} \left\{ \sum_{i=1}^{n-1} V_i(\hat{\mathbf{r}}) \frac{\partial E}{\partial x_i} + \frac{\partial E}{\partial t} \right\}. \quad (14)$$

For this case we propose the following strategy (again we define the curve as an integral curve of the vector field V_i):

$$(a) (\forall i) \frac{dx_i}{dt} = H(t), \quad (b) H(t) \sum_{i=1}^{n-1} \frac{\partial E}{\partial x_i} + \frac{\partial E}{\partial t} = H(t). \quad (15)$$

Hence the curve has the form $x_i = \int H(t) dt$, ($i = 1, \dots, n-1$). The function E is determined by a first order partial differential equation of a certain special form (Eq.(15b)).

One may think that the way in which we have solved the problems is artificial because we introduced *ad hoc* vector fields in the reasoning. This is not really the case, it is just the effect of our rigid vision of the process of solution.

Consider, for example, the Poincaré-Cartan *1-form* of classical mechanics:

$$W = \sum_{i=1}^n p_i dq_i - H dt.$$

We do not have the right to write down it as $W = dS$, where S is the action, until we prove that it is in fact an integrable *1-form*. With this purpose in mind we can attack the problem in the following way: we suppose that the form is integrable and we write $p_i = \partial S / \partial q_i$, $H(p_i, q_i) = -\partial S / \partial t$, and we get the Hamilton-Jacobi equation. Hence, the problem of integrability is the problem of the existence of solutions of the Hamilton-Jacobi equation. As it is well-known, an analytic solution for this equation always exists locally (Cauchy-Kovalevsky theorem), hence, the *1-form* is a locally integrable *1-form*. In the dynamical problem we know the Hamiltonian explicitly; but we know *neither* the form of the curve *nor* the action as a function of the coordinates (not as a functional, because that is another point of view). But, as it is well-known, if we can solve the Hamilton-Jacobi equation we know the action and the solution of the dynamical problem by means of a canonical transformation generated by this action function. Clearly, in this case we have introduced all our “auxiliary functions”, the action and the Hamiltonian, to know the explicit form of the curve in phase-space. Of course, we have required some data: the form of the Hamiltonian and the supposition of integrability of the *1-form*. And from the

theoretical point of view, it is enough to construct the solution of the dynamical problem. However, we need to make our distinction in this point: *the action as a function of the coordinates differ from the function constructed by restriction of the action to the curve.*

Another important point becomes clear when we use *1-forms*: in all the cases which we have treated, we need to suppose the integrability of a *1-form*. For example, when we treat the $\{E\}$ -case we start from the *1-form*:

$$dE = \sum_i \frac{\partial E}{\partial x_i} dx_i + \frac{\partial E}{\partial t} dt$$

which is clearly integrable. Hence, we want to know a curve as an integral curve of a vector field which we define as:

$$X = \sum_{i,j} b_{ij} \frac{\partial E}{\partial x_j} \frac{\partial}{\partial x_i} + \frac{\partial}{\partial t}.$$

The inner product of these two tensors (the pairing between the tangent and cotangent space) give us the result:

$$\langle dE, X \rangle = \frac{\partial E}{\partial t}(x_1, \dots, t),$$

hence, the composition is in fact, the result of taking the limit of the inner product in the integral curves of the vector field X . We can treat the other cases from this point of view, but that is easy after this explanation. In a geometric interpretation we have the following elements: the tangent vectors, and the angle between them. In the $\{E\}$ -case we have the normal, but we have neither the tangent nor the angle; in the $\{\mathbf{p}\}$ -case the tangent, but we have neither the normal nor the angle; finally, in the $\{f\}$ -case we have the angle, but we have neither the tangent nor the normal.

Now let us make a brief review of the last case ($\{\}$ -case).

The point is that in this case we have no any data and to treat it we need some information. Heyting [3] notes that we ought to distinguish two different concepts:

1. Theories of the constructible.
2. Constructive theories

The first one is characterized by 3 conditions:

- (a) we presuppose a mathematical theory in which the class of constructible objects can be defined;

- (b) the notion of a constructibility is no primitive;
- (c) we have a liberty to choose the definition of a constructible, But, of course, it must correspond to our intuitive notion of a mathematical construction.

For the second point (the constructive theories) Heyting says: “*a theory in which an object is only considered as existing after it has been constructed. In other words, in a constructive theory there can be no mentioning of other than constructible objects*”. The main feeling of Heyting is expressed in the following sentence: “*I am unable to give an intelligible sense to the assertion that a mathematical object which has not been constructed exists.*”

In the case which we want to treat we have no any data concerning the equation $f = E \circ \mathbf{p}$. Hence, if we accept that we can only speak about those objects which can be constructed explicitly (or, at least, we have a method to construct them), the case which we are treating, the $\{\}$ -case, is not even a case. It is nothing, it is just a line of symbols without any meaning. For this reason when one speaks about the functional equation $f = E \circ \mathbf{p}$ one, in fact, is speaking about the cases considered before: $\{E, \mathbf{p}\}$ -, $\{E, f\}$ -, $\{\mathbf{p}, f\}$ -, $\{E\}$ -, $\{f\}$ -, and $\{\mathbf{p}\}$ -case.

As a last remark we can see that we have shown that the *generally accepted* expressions of the type of Eq.(3) *cannot be valid*.

III. ABOUT FUNCTIONAL EXTENSIONS

We shall give our problem the most general setting. Let us start with a topological space D , so that it is possible to construct the general object of arrows: $T(D, K)$ where T is any covariant functor. Hence we can construct the functor:

$$T(D, *) : \mathbf{C}_1 \rightarrow \mathbf{C}_2, \quad (16)$$

where $\mathbf{C}_i (i = 1, 2)$ are any small categories. Then for each arrow we have $f \in T(D, K)$ the diagram: $f : D \rightarrow K$. For us, the following situation is the most important: the set D is an object with a given structure, so we use the symbol $P(D)$ to denote its power set (which is a topology, of course, any topology is a subset of the power set, but not any subset of the power set is a topology). In this way, for each element in $P(D)$ we can define the following elements: $\langle f_A, A \rangle$ for all $A \in (P(D))$. Here the symbol $\langle f_A, A \rangle$ means

that the object $A \in P(D)$ is put in correspondence with the function f_A . So we may form the set:

$$F_D = \{\langle f_A, A \rangle | A \in P(D)\} \quad (17)$$

of functional elements.

It is clear that this procedure has been realized in a somewhat formal manner, however, this is the more general form. As we can see, there are several elements which are important for our construction: the covariant functor T , the object D , its power set $P(D)$, the set of elements F_D which is the part in which the functions enter the discussion and the two small categories: $\mathbf{C}_1, \mathbf{C}_2$.

Definition 1: We shall call the symbol $\langle F_D, T, \mathbf{C}_1, \mathbf{C}_2, P(D) \rangle$ a general function.

The idea behind a general function is that all its elements are different for each element of the power set of D . Sometimes we can use a specific topology instead of the power set, but this choice relies on our convenience. Besides, we can see that in general, any topology is just a subset of $P(D)$. The formation of a topology in the object D can respect its structure or not. We use the categorical notions to introduce the generality which they carry, because in general a function depends on the categories in which it is defined, see [4] chap. 1, for more details. Hence, the general setting is: *how are the different elements of a general function related?* The problem may seem trivial without more elaboration, however, as we have seen in the introduction, in some realms the problem is not trivial. Let us give a few additional examples.

Example 1: Consider the following example (see [5]), which is clearly not trivial, suppose the following choice: $D = \mathbf{C}$ where \mathbf{C} is the complex plane, if we use the symbol **Anal** to denote the functor of the set of complex analytic functions we have:

$$\mathbf{Anal}(\mathbf{C}, *) : \mathbf{Set} \rightarrow \mathbf{Set} \quad (18)$$

or, to be more concrete, the arrows: $f : \mathbf{C} \rightarrow \mathbf{C}$ of complex analytic functions are at hand. We must consider the power set $P(\mathbf{C})$ of the complex plane and the construction of the elements: $\langle f_A, A \rangle$ for each element of the power set. This is the most general situation for the choice that we have made of our basic elements. In this setting we have the following group of well-known definitions [5]:

Given two “functional elements” $f(A) = \underset{\text{def}}{\langle} f_A, A \rangle$, $f(B) = \underset{\text{def}}{\langle} f_B, B \rangle$, we can say that we have a direct analytic prolongation if, and only if, the following two conditions hold:

$$A \cap B \neq \emptyset$$

$$f_A = f_B \quad \text{in } A \cap B$$

So, we can see that in general, the problem of analytic continuation is a realization, in the complex domain, of our definition of a general function

Example 2: Consider the functors:

$$C(D, *) : \mathbf{Top} \rightarrow \mathbf{CRng} \quad \text{and} \quad C^*(D, *) : \mathbf{Top} \rightarrow \mathbf{CRng} \quad (19)$$

from the topological spaces to the rings of continuous functions. The functor C^* is for bounded functions. Here the problem is as follows: a set S is C -embedded if, and only if, every function $f \in C(S)$ can be extended to a function $g \in C(D)$. Here $S \subset D$ and $C(S)$ is an abbreviation of $C(S, S)$. The idea here is that the extension is a C -function. The definition of C^* -embedding is similar.

One of the most important characteristics of the C^* -embedding is Uryshon’s **lemma**:

A subset S of the set D is C^* -embedded in D if, and only if, any two completely separated sets in S are completely separated in D .

We can see that this lemma is just an assertion about functional extensions, that is, a theorem about the way in which the elements of a general function are related [6], p. 18. Here the set F_D can be constructed once we have fixed the topology of the spaces S and D , or at least the base of the topology. If we use the power set we have a conceptual generality, but we can fall into troubles for some purposes. Let us take for topology of D its power set $t(D)$, hence, the set F_D can be formed and we have:

$$F = \langle F_D, C, \mathbf{Top}, \mathbf{CRng}, t(D) \rangle \quad (20)$$

as our general function for this case. Of course, we can construct F without recourse to the Uryshon’s lemma, however, this result gives us a way to relate two elements of the general function F .

Example 3: Let us come back to the example in the introduction. Consider the functor:

$$C^\infty(R^n, *) : \mathbf{Vect} \rightarrow \mathbf{Vect} \quad (21)$$

so the arrows: $f : R^n \rightarrow R^n$ where R is the real line. The power set is now $P(R^n)$, and the set of functional elements is: $\{\langle f_A, A \rangle | A \in P(R^n)\}$. The notion of differentiability does not change and we can define the derivative of a general function as the general function formed with the derivatives of the functional elements of the starting general function. If one of such elements is not differentiable, the general function is not.

IV. THE LIMITING PROCEDURE

Let us explain in more detail the limiting procedure which can be used for the elements of a general function. Consider the initial object D and suppose a partition of the form:

$$D = \bigcup_{i=1}^n G_i \quad (22)$$

Hence, the general function is defined with the help of the elements of the set $F_D = \{\langle f_i, G_i \rangle\}$ and the functor $T(D, *)$. We can make this decomposition in many ways. For example: $D = \bigcup_{A \in t(D)} A$, where $t(D)$ is the power set of D .

Now we define the system of sets:

$$P_i = \{A \in P(D) | (G_i \subset A)\}. \quad (23)$$

In other words: the set of all the sets A so that the set G_i is contained. Is very easy to show that each P_i is a filter.

Lemma: Each P_i is a model of a filter in D (see [6] p. 24).

Proof: \diamond (a) We can see that \emptyset is not an element of P_i , any i , because if $\emptyset \in P_i$ then we can find a set A such that $A \subset \emptyset$ which is a contradiction, hence we have proved the first axiom. (b) If we suppose that $A, B \in P_i$ then $A \cap B \in P_i$ because, at least A and B have the set G_i in common, hence G_i is in their intersection, but this is the condition for belonging to P_i . The second axiom is satisfied. (c) If $A \in P_i$, $B \subset D$ and $A \subset B$ it is very easy to see that $B \in P_i$. The lemma is proved. \diamond

This lemma (such trivial as it is) is important, because with a filter we can define a limit for the elements of a general function. In fact, given the filter P_i of the element i

in the partition, we have a function f_i which maps the element G_i . Clearly, the elements $f_i(G_i)$ are the images of the set D by the general function F_D . So, we define the filter of the image of D by the general function F_D as $F_D(P_i) = \{A \in K | f_i(G_i) \subset A\}$. This is clearly a filter. With these elements it is possible to set up a well-known definition of the limiting procedure for the elements of a general function.

Definition 2: A set G_i is the limit of a filter H if, and only if, H is *stronger* than the filter P_i .

Definition 3: Consider the filter P_i , hence the set $A \subset K$ is the limit of the general function F_D under the filter P_i if, and only if the set A is the limit of the filter $F_D(P_i)$. That is, without abbreviations: the filter $F_D(P_i)$ is stronger than the filter formed with the sets that contain A .

We say that a filter H is stronger than the filter B (of course, both filters defined on the same space) if, and only if, for any $a \in B$ there is a set $b \in H$ such that $b \subset a$. Of course, this is just the notion of approximation, because a filter H is stronger than a filter B if their elements are nearest to a certain set than the elements of B . Now let us use this concept for the example in the introduction. We have the equation:

$$\frac{d}{dt} E_\gamma = \lim_{\mathbf{r} \rightarrow \mathbf{r}(t)} \left\{ (\mathbf{V} \cdot \nabla) E_A + \frac{\partial}{\partial t} E_A \right\}, \quad (24)$$

where E_γ is a function along the curve γ and E_A is a function defined on the set A . Now, let us give a precise meaning to the process involved. We have the general function E_D and two of their functional elements are involved: $\langle E_\gamma, \gamma \rangle$, $\langle E_A, A \rangle$ where γ and A are sets in D . Hence we can see that the limiting procedure affects only the functional element $\langle (\mathbf{V} \cdot \nabla) E_A, A \rangle$ so, we do the following: we select a set $d \subset \gamma$ and we form its filter P_d , so, a set $B \in R^n$ in the image of the functional element $(\mathbf{V} \cdot \nabla) E_A$ is its limit if, and only if the filter formed with the sets that contain the image of $(\mathbf{V} \cdot \nabla) E_A$ is stronger than the filter formed with the sets that contain B . Of course the extension of this definition covers the usual ε - δ arguments.

V. CONCLUSIONS

As promised in the introduction, we have solved in its most general form the ambiguity which arises in the differential calculus of several variables with the help of category theory. Besides we have showed several examples of realizations of our construction.

ACKNOWLEDGMENTS

The authors would like to express their gratitude to Prof. Valeri Dvoeglazov for his discussions and critical comments. We would also like to thank Annamaria D'Amore for revising the manuscript.

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